

(ALGEBRA)

511.4
036

LECTURES DELIVERED TO POST-GRADUATE STUDENTS OF
CALCUTTA UNIVERSITY

BY

FRIEDRICH WILHELM (LEVI) DR. PHIL. NAT.
HARDENCE PROFESSOR

PART I

SYSTEMS OF LINEAR EQUATIONS



PUBLISHED BY THE
UNIVERSITY OF CALCUTTA
1936

BCU 1738

PRINTED AND PUBLISHED BY SHYAMPRASAD SARKAR
AT THE CALCUTTA UNIVERSITY PRESS, SENATE HOUSE, CALCUTTA.

Vol. No. 988B—August, 1936—1

102.330



PREFACE

On introducing a new course of lectures in Algebra I realized after delivering a few lectures that the students of this country should have in their hands a book covering the whole subject-matter of the lectures. In order to make my lectures successful I had no other alternative than to write a text-book and to publish it in different parts as quickly as possible.

So a provisory edition of this text-book is taken in hand. References to the original papers, examples, explanation of details, everything that causes delay of publication had to be omitted in these "lectures." Later on a full text-book on Algebra will be published.

In placing this fascicule in the hands of the students, I offer my heartiest thanks to our energetic Vice-Chancellor, SYAMAPRASAD MOOKERJEE, Esq., M.A., B.L., Barrister-at-Law, M.L.C., without whose sympathetic co-operation this publication would not have come into being. I thank also the Calcutta University Press for having printed this paper in a very short time under difficult circumstances.

Proofs and manuscripts have been revised by Mr. R. C. Bose, M.A., Mr. S. K. Bhar, M.Sc., and especially by Mr. A. C. Choudhury, M.Sc. If the reader do not find many offences against the spirit of the English language, he should be thankful to these three young colleagues of mine.

CALCUTTA, ASUTOSH BUILDING.

P. W. LEVI.

August, 1936.

Remark: Division is permissible, if and only if the denominator is not equal to zero. If the denominator is not a constant, but a function, it is necessary to treat separately the cases in which this function vanishes. In the lecture, examples of this kind will be given.

II. $n=2, m=1$ $a_1 x_1 + a_2 x_2 = a_0$

(a) $a_1 \neq 0, a_2 \neq 0$. Solutions (x_1, x_2) : one of the numbers is arbitrary, the other is defined by it.

(b) $a_1 \neq 0, a_2 = 0$. Solutions $(a_0 : a_1, x_2)$ x_2 is arbitrary.

(c) $a_1 = 0, a_2 \neq 0$. „ $(x_1, a_0 : a_2)$ x_1 „

(d) $a_1 = a_2 = 0, a_0 \neq 0$. No solution.

(e) $a_1 = a_2 = a_0 = 0$. Solutions x_1 and x_2 are arbitrary.

III. $n=1, m=2$ $a_1 x = a_0$

$$b_1 x = b_0$$

(a) $a_1 b_0 - b_1 a_0 \neq 0$. No solution.

(b) $a_1 b_0 - b_1 a_0 = 0, (a_1, b_1) \neq (0, 0)$.

Solutions: $a_0 : a_1$, or $b_0 : b_1$, or $a_0 : a_1 = b_0 : b_1$

if $a_1 \neq 0$ $b_1 \neq 0$ $a_1 \neq 0$ $b_1 \neq 0$.

(c) $a_1 = b_1 = 0, (a_0, b_0) \neq (0, 0)$. No solution.

(d) $a_1 = b_1 = a_0 = b_0 = 0$. Solution x is arbitrary.

IV. $n=2, m=2$ $a_1 x_1 + a_2 x_2 = a_0$

$$b_1 x_1 + b_2 x_2 = b_0$$

$$a_1 b_2 - a_2 b_1 = \Delta$$

$$a_2 b_2 - a_2 b_0 = \Delta_1$$

$$-a_0 b_2 + a_1 b_0 = \Delta_2$$

Necessary conditions for solutions (x_1, x_2) are $\Delta x_1 = \Delta_1$
 $\Delta x_2 = \Delta_2$

- (a) $\Delta \neq 0$. One solution $(\Delta_1 : \Delta, \Delta_2 : \Delta)$
 (b) $\Delta = 0$. $(\Delta_1, \Delta_2) \neq (0, 0)$. No solution.
 (c) $\Delta = \Delta_1 = \Delta_2 = 0$. $a, b, (=b, a) \neq 0$.

Every solution of the first equation is also a solution of the second.—
 See II (a).

- (b) $a_1 = a_2 = 0, a_3 \neq 0$. No solution.
- (b') $b_1 = b_2 = 0, b_3 \neq 0$. No solution.
- (c) $a_1 = a_2 = a_3 = 0$. Every solution of the second equation is a solution of the system.—See II.
- (c') $b_1 = b_2 = b_3 = 0$. Every solution of the first equation is a solution of the system.—See II.
- (d) $a_1 = b_1 = \Delta_1 = 0$. Solutions.—See III, β, γ, δ .
- (d') $a_2 = b_2 = \Delta_2 = 0$. Solutions.—See III, β, γ, δ .

§ 2. THE HOMOGENEOUS SYSTEM BELONGING TO AN ARBITRARY SYSTEM OF LINEAR EQUATIONS.

From the examples given in § 1 we may realize that even in the most simple cases a great number of different sub-cases has to be considered. The number of these sub-cases seems to increase infinitely with n and m . For avoiding these difficulties we will follow a way very characteristic of mathematical thoughts. We will suppose that solutions have been found out, and we will enquire about the connection between those solutions. By considering the properties of the system of all solutions we will get a general theory of the systems of linear equations including also different ways for finding out the solutions of a given system.

If $(x_1, x_2, \dots, x_n) = (w_1, w_2, \dots, w_n)$ and $(x_1, x_2, \dots, x_n) = (u_1, u_2, \dots, u_n)$ are solutions of (2), then $(x_1, x_2, \dots, x_n) = (w_1 - u_1, w_2 - u_2, \dots, w_n - u_n)$ is a solution of

$$k_1 x_1 + \dots + k_n x_n = 0$$

and conversely, if (u_1, \dots, u_n) is a solution of (2), and (y_1, \dots, y_n) a solution of (2/H), then $(u_1 + y_1, \dots, u_n + y_n)$ is a solution of (2). Therefore the following theorem holds:

Theorem 1. Starting from an arbitrary solution of (2) we will get all solutions by addition of the solutions of (2/H).

The homogeneous system (2/H) belongs to the system (2). For solving (2) we have to find out all solutions of (2/H) and an arbitrary one of (2). For that purpose the introduction of a new notion is convenient.

§ 3. THE n -vectors

Definition 1: An ordered set of n numbers is called an n -vector.

$$\alpha = (\alpha_1, \dots, \alpha_n) \quad (3)$$

The n numbers α_i defining the n -vector are called its co-ordinates. As this set is an ordered one, the vector will generally be changed by the interchange of the co-ordinates.

Examples: (1) The co-efficients of an arbitrary equation of (2/H) define an n -vector; it is called the "vector of that equation," and also the "vector of that row."

(2) The co-efficients of an arbitrary column of (2/H) define an m -vector, the "vector of that column."

(3) The solution (1) of (2) defines an n -vector, the "vector of the solution."

(4) Vectors in the plane (the space), in the sense this word is ordinarily used, are 2-vectors (3-vectors).

(5) Let n be the number of the customers of a bank; the balances of the customers are the co-ordinates of an n -vector representing the actual state of the bank. The reader may interpret the vector-addition defined below for this example.

Definition 2: The product of a number c and the vector α is an n -vector.

$$c\alpha = (c\alpha_1, \dots, c\alpha_n). \quad (4)$$

Definition 3: The sum of α and $\beta = (\beta_1, \dots, \beta_n)$ is

$$\alpha + \beta = (\alpha_1 + \beta_1, \dots, \alpha_n + \beta_n). \quad (5)$$

From these definitions it follows:

$$\begin{aligned}
 \alpha + \beta &= \beta + \alpha && \text{commutative law,} \\
 \alpha + (\beta + \gamma) &= (\alpha + \beta) + \gamma && \text{associative law,} \\
 c(\alpha + \beta) &= c\alpha + c\beta && 1^{st} \text{ distributive law,} \\
 (c_1 + c_2)\alpha &= c_1\alpha + c_2\alpha && 2^{nd} \text{ distributive law,}
 \end{aligned} \tag{6}$$

As these laws hold, we can use the notations of sums of n -vectors in the same manner as it is to be used for numbers:

$$\sum c_i \alpha^i = (\sum c_i \alpha_{1i}^1, \dots, \sum c_i \alpha_{ni}^n) \tag{7}$$

$$i = 1, \dots, m; \text{ } c \text{ being arbitrary numbers; } \alpha^i = (\alpha_{1i}^1, \dots, \alpha_{ni}^n)$$

being arbitrary n -vectors.

The vector $-1 \cdot \alpha$ is called the *negative* of α and written $-\alpha$. (The addition of $-\alpha$ is the inverse operation to the addition of α . As in elementary arithmetics this inverse operation is called *subtraction* and written by the sign $-$. Therefore:

$$\beta + (-\alpha) = \beta - \alpha. \tag{8}$$

The following special n -vectors will often be used:

$$\begin{aligned}
 0 &= (0, \dots, 0) && \text{Zero-vector} \\
 e^1 &= (1, 0, \dots, 0) && 1^{st} \text{ Unit-vector,} \\
 e^2 &= (0, 1, 0, \dots, 0) && 2^{nd} \text{ Unit-vector} \\
 &\dots\dots\dots && \\
 e^n &= (0, 0, \dots, 0, 1) && n^{th} \text{ Unit-vector.}
 \end{aligned} \tag{9}$$

$$\text{Formulae: } \alpha - \alpha = 0$$

$$c \cdot 0 = 0 \tag{10}$$

$$\alpha = \sum a_i e^i$$

§ 4. VECTOR-SPACES.

Definition 3: The n -vector $\sum c_i \alpha^i$ is dependent on the n -vectors α^i .

Definition 4: The n -vectors $\alpha^1, \dots, \alpha^m$ are said to be independent if none of them is dependent on the $m-1$ others and $m > 1$. A single n -vector is independent, if it is $\neq 0$.

Definition 5: The set of all vectors dependent on the α^i is called the vector-space generated by the α^i .

Definition 6: A set of independent n -vectors generating a vector-space is called its Basis.

Theorems concerning vector-spaces:

✓ 1. The n -vectors α^i are independent if and only if for every system of numbers $(c_1, \dots, c_m) \neq (0, \dots, 0)$ we have $\sum c_i \alpha^i \neq 0$.

Proof. 1. If $c_i \neq 0$, α^i is dependent on the other α^j . 2. If α^i is dependent on the other α^j , then $\alpha^i = \sum d_j \alpha^j$, and therefore $\sum c_i \alpha^i = 0$, for $c_i = -1$, $d_j = c_j$. The theorem is evident in the case of a single n -vector.

✓ 2. If β^1, \dots, β^r belong to a vector-space V , every n -vector α dependent on the β^i belongs also to V . *and they do not form the basis*

Proof. Let V be generated by $\alpha^1, \dots, \alpha^m$, then from $\alpha = \sum k_i \beta^i$, $\beta^i = \sum b_j \alpha^j$ follows $\alpha = \sum c_j \alpha^j$, where $c_j = \sum k_i b_j$.

✓ 3. Every vector-space containing n -vectors $\neq 0$ has a basis.

Proof. If the vectors α^i generating V are not independent, α^i may depend on the other $m-1$ generating n -vectors. From 2 it follows that these vectors generate also V . On repeating—if it is necessary—this reduction we will get after a finite number of steps a subset of the α^i generating V composed of independent vectors, i.e., a basis of V .

✓ 4. If $\alpha^1, \dots, \alpha^m$ is a basis of V , $\beta = \sum c_i \alpha^i$, and $c_j \neq 0$, then we will get a new basis of V on replacing α^j by β .

Proof. The vector-space generated by β and the $\alpha^{i \neq j}$ is contained in V . On the other hand it contains $\beta - \sum_{i \neq j} c_i \alpha^i = c_j \alpha^j$ and therefore α^j .

We have to show that the m n -vectors β and $\alpha^{i \neq j}$ are independent.

Let $d\beta + \sum_{i=1}^m d_i \alpha^i = 0$, on replacing β by its value $\sum_{i=1}^m c_i \alpha^i$ we get a vanishing linear function of the α^i whose coefficients vanish, as the α^i are independent. The coefficient of α^j is $d c_j$; as $c_j \neq 0$, it follows: $d=0$. As the α^i are independent, the d_i are vanishing. Therefore the α_i and β are independent; so they form a basis of V .

5. If $\alpha^1, \dots, \alpha^m$ is a basis of V , and β^1, \dots, β^r are independent vectors in V , then we get a new basis on replacing t suitable α^i by the β , and $t \leq m$ holds.

Proof. From 4 it follows that a suitable α can be replaced by β^1 . Let $\alpha^1, \dots, \alpha^{t-1}, \beta^1, \dots, \beta^r$ be a basis of V , $r < t$ then we can express β^{r+1} by these m vectors, and in this expression the coefficient of at least one α do not vanish; therefore this α can be replaced by β^{r+1} . Therefore we can continue replacing an α by a β till $r=t$, or $r=m$. In the last case all β 's must depend on the β^1, \dots, β^r , and therefore $t=m$ in this case. Generally $t \leq m$.

Definition 7: The maximum number of independent vectors of V is called the Rank of V .

6. The number of the vectors of an arbitrary basis of V equals the rank of V . (Therefore every basis of V has the same number of vectors.)

Proof. From 5 it follows that there cannot be more independent n -vectors in V than an arbitrary basis has elements.

7. If every n -vector of V belongs to V' , but not every n -vector of V' belongs to V , then is the rank of V less than the rank of V' .

Proof. Let $\alpha^1, \dots, \alpha^r$ be a basis of V , and β a vector of V' not contained in V , then is r the rank of V , but the rank of V' is at least $r+1$ because V' contains $r+1$ independent elements α^1, β .

8. The rank of a vector-space of n -vectors is at most n .

Proof. The n unit-vectors (9) generate a vector-space V' in which the co-ordinates are arbitrary numbers; therefore V' contains every n -vector, and 8 follows from 7.

9. Between $p > n$ of n -vectors there exists always a linear equation with non-vanishing coefficients.

Proof. If these n vectors are independent they will generate an n -vector-space of rank $p \geq n$.

1. A system A of n vectors with the property that the sum of two arbitrary elements of A , as well as the product of an arbitrary element of A with an arbitrary real number belongs to A is a vector space.

Proof. From it follows that in A there exists a maximum number r of independent vectors. The vector space generated by such r independent vectors is identical with A . If $r=0$ A contains only the n -vector 0 .

4. THE VECTOR SPACE ASSOCIATED WITH A SYSTEM OF HOMOGENEOUS LINEAR EQUATIONS

THEOREM II. The solutions of $2/H$ form a vector space X .

Proof. From $\sum_{j=1}^n a_{ij}x_j = \sum_{j=1}^n b_{ij}x_j = \dots = \sum_{j=1}^n k_{ij}x_j = 0$

and $\sum_{j=1}^n a_{ij}x_j = \sum_{j=1}^n b_{ij}x_j = \dots = \sum_{j=1}^n k_{ij}x_j = 0$ it follows

$$\sum a_{ij}cx_j = \sum b_{ij}cx_j = \dots = \sum k_{ij}cx_j = 0 \quad \text{and}$$

$$\sum a_{ij}(x_j + y_j) = \sum b_{ij}(x_j + y_j) = \dots = \sum k_{ij}(x_j + y_j) = 0$$

Therefore the solutions satisfy the conditions of a vector space given in §4, 10.

THEOREM III. In every vector space X there exists a vector space V such that every n -vector of X is a solution of a linear homogeneous equation f and only if the vector of this equation is a vector of V .

Proof. The vectors of X are solutions of the equation defined by the n -vector α . If they are solutions of the equations defined by α and by β then they are also solutions of the equations defined by $\alpha + \beta$ and by $\gamma\alpha$. From §4, 10 it follows that the set of the vectors α, β, \dots form a vector space.

The Theorems II and III show that our problem of II is closely connected with two vector spaces X and V . The vectors of $2/H$ generate a vector space V' and every vector of V' is also a vector of V . If V and V' were not identical then V should be of higher rank. In that case two vector spaces of different ranks would define the same

vector space V . In other words there would exist a vector λ independent of the vectors of the equations (2.11), such that every solution of (2.11) is also a solution of the equation generated by λ . We will see later that this case is impossible.

10. THE BASIS OF A VECTOR SPACE - THE METHOD OF SUBSTITUTION

For further consideration of linear equations we need know some more properties of vector spaces.

11. If $\alpha, \beta, \dots, \kappa$ are n vectors generating a vector space V and a, b are arbitrary numbers then $a\alpha + b\beta, \dots, \kappa$ generate also V .

Proof. As V includes $a\alpha + b\beta$ and the vector-space generated by $a\alpha + b\beta, \beta, \dots, \kappa$ includes α the two vector spaces are identical.

12. If the n vectors $\alpha, \beta, \dots, \kappa$ generating V are not all equal to 0, then we will get another system of n vectors generating V by omitting the m vectors equal to 0.

Proof. The m vector 0 is dependent on the other ones.

13. The n -vectors

$$\begin{aligned} \beta^1 &= (1, 0, \dots, 0, 0, b^1_1, \dots, b^1_n) \\ \beta^2 &= (0, 1, \dots, 0, 0, b^2_1, \dots, b^2_n) \\ &\vdots \\ \beta^m &= (0, 0, \dots, 0, 1, \dots, b^m_1, \dots, b^m_n) \end{aligned} \quad (11)$$

are independent

Proof. Let $\sum \beta^i = 0$ the coordinates of the vector on the left are equal to the coordinates on the right hand. Therefore $b^i_1 = \dots = b^i_n = 0$.

14. V may be an arbitrary vector space including n vectors $\neq 0$, starting from an arbitrary finite system of n vectors generating V and using the methods of 11 and 12 we will get after a finite number of suitable steps a basis of V differing from the n vectors of formula 11 at most by a permutation of the coordinates of all vectors.

Remark. The methods 11 and 12 can be considered as operations practised on the scheme of coordinates written on the right hand in formula

11. It is therefore useful to take (11) as a property of that scheme. For this purpose we introduce a notation often used in mathematics.

Definition 8. The system of the coordinates of an ordered n -vector is called a M -matrix M . The coordinates of the i -th vector form the i -th row of M , the i -th coordinate of the n -vector forms the i -th column of M ($i = 1, \dots, n$). $M=0$ means that every co-ordinate of M vanishes.

The operations given by 11 and 12 may be shortly expressed by the words "row addition" and "omission of rows". Using these words we get a theorem (14) equivalent to (14').

14. By the operations of row addition and omission of 0 rows every matrix $M \neq 0$ can be transformed to a matrix of the kind (14) or to a matrix differing from (14) by a permutation of the columns. The proof of 14' and therefore also the proof of 14 can be made by different steps. To simplify the description of these transformations it may be understood that every row that would appear should be omitted automatically without mentioning this operation. Every row contains therefore at least one co-ordinate $\neq 0$ for all transformations.

the rows of the matrix (1), (2)

the coordinates of (1) $\{1, \dots, r, \dots, n\}$, $r = 1, \dots$ (12)

the columns $\{1\}, \{2\}, \dots, \{n\}$

These signs do not denote constant rows. They change at every step. If r is the first row $r = 1$ denotes the first row at every step, also $\{r\}$ the r -th column, n being an arbitrary number, etc.

Choosing these signs the different steps of the transformation may be described as follows:

1. a) r may be the same or a number such that $\{r, 1\} \neq 0$

by the row addition $(1) \rightarrow (-1) \cdot \{r, 1\}$ (1)

$\{r, 1\}$ becomes -1

1. b) by the row addition $(k) \rightarrow k \cdot \{r, k\} - 1$

$\{r, k\}$ becomes 0

By these permutations $\langle i_1 \rangle$ and $\langle i_2 \rangle$ have been swept out, $\langle i_1 \rangle$ one of the numbers has been made to be zero or vanish. We continue on sweeping out.

2. a_{12} may be the smallest number such that $a_{12} \neq 0$ then in $\langle i_2 \rangle$, by row addition $2) \rightarrow \{-1, \dots, a_{12}, \dots\}$ (2)

$[i_2, 2]$ becomes -1 and $[i_2, 2] = 0$ is obtained.

3. By the row addition $1) \rightarrow \dots + [i_2, 2] \cdot 2$

for every $k \neq 2$, $[i_2, k]$ becomes 0 and $\langle i_2 \rangle$ is not changed.

Now the columns $\langle i_1 \rangle$ and $\langle i_2 \rangle$ are swept out. We continue the sequence of operations. After $2(i_1 - 1)$ steps the columns

$\langle i_1 \rangle, \langle i_2 \rangle, \dots, \langle i_{r-1} \rangle$ may be swept out. i.e.

$[i, s] = -1$ $[i, k] = 0$ for $k = 1, \dots, q-1$ if $s = i_1, \dots, i_{r-1}$

being different numbers. If after these steps (and the elimination of 0 rows) the matrix has more than $r-1$ rows

4. a_{1r} may be the smallest number for which $[i_r, 1] \neq 0$

by $(1) \rightarrow \{-1, [1, q], \dots, [i_r, q]\}$ becomes 1

5. By $(1) \rightarrow 1) + [i_r, q] \cdot [i_r, 1]$ becomes

the rows $\langle i_1 \rangle, \dots, \langle i_{r-1} \rangle$ are not changed by these transformations.

By mathematical induction it follows that the method can be continued till the number of columns is equal to the number of the remaining rows. After this transformation

$\langle i_1 \rangle, \dots, \langle i_r \rangle$ may be swept out $[i, s] = -1$ if $s = i_1, \dots, i_r$

for $s = 1, \dots, r-i+1$. By a permutation of the columns transforming $i_r \rightarrow 1$ the matrix will take the form 1).

Conclusion. Starting from an arbitrary system of n -vectors generating $V \neq 0$ it is always possible to get a basis of V .

Proof. Let m be the number of the generating vectors. It is not possible to sweep out the matrix by at most m^2 row additions and omissions of at most $m-1$ rows. As the n -vectors of the remaining rows are independent, they form a basis.

6.7. SOLUTION OF SYSTEMS OF HOMOGENEOUS LINEAR EQUATIONS BY THE METHOD OF "SWEEP-OUT"

Theorem IV. Let V be the vector space generated by the n vectors of H , X the vector space of the solutions of H , and V' the vector space defined by (13).

$$\text{Then } V = V', \text{ rank } (V) = \text{rank } (X) = n - \text{rank } (H).$$

Theorem V. It is possible to find out a basis of X by a limited number of row additions and omissions of 0 rows.

Proof of the Theorems IV and V. Every solution of H is a solution of the equations belonging to the n vectors of V and conversely. To obtain these solutions it is sufficient to find out the solutions of the equations belonging to a basis of V . As it was proved by 14 it is always possible to find out such a basis by the method of sweep-out. This basis can be taken in the form (1) only by a permutation on $i_1 \rightarrow k$ of the columns. If we write therefore $x_i = y_k$, $k=1, \dots, n$

then our problem will be reduced to the following

$$\begin{aligned} y_1 + b_{12} y_2 + \dots + b_{1n} y_n &= 0 \\ y_2 + b_{22} y_2 + \dots + b_{2n} y_n &= 0 \end{aligned} \quad (14)$$

$$y_i = \sum_{j=1}^{n-r} b_{ij} y_j, \quad i=1, \dots, r$$

The y_1, \dots, y_r can take arbitrary independent values and y_{r+1}, \dots, y_n are completely determined by them. There exists therefore a (1, 1) correspondence between the vector space of $n-r$ vectors (y_{r+1}, \dots, y_n) , and the vector space of all solutions (x_1, \dots, x_n) .

$$\begin{aligned} (1, 0, \dots, 0, 0) & \quad (b_{12}, \dots, b_{1n}, 1, 0, \dots, 0, 0) \\ & \quad \text{and} \\ (0, 0, \dots, 0, 1) & \quad (b_{22}, \dots, b_{2n}, 0, 0, \dots, 0, 1) \end{aligned}$$

are corresponding vectors. As this correspondence is invariant for vector addition and multiplication with a number, and the $n-r$ vectors on the right form a basis of the solutions (14), we will get a basis of X by the permutation $i \rightarrow k$ of the coordinates.

Therefore Theorem V is true, and

$$\text{Rank } (V) + \text{rank } (X) = n$$

leads. As X is also the vector space of all solutions of V we can replace V' by V in this formula, hence (14) holds. Every n vector of V' belongs to V and these vector spaces have the same rank. Therefore (14, 2) V and V' are identical.

§ 8. SOLUTIONS OF NON-HOMOGENEOUS LINEAR SYSTEMS

In order to solve (2) let π and l be $(n+1)$ and $(2n+1)$, and

$$\begin{aligned} l_1 x_1 + \dots + l_n x_n + l_{n+1} x_{n+1} &= 0 \\ l_1 x_1 + \dots + l_n x_n + l_{n+1} x_{n+1} &= 0 \end{aligned} \quad (15)$$

The generating $(n+1)$ vectors being called

$$\alpha = (\alpha_1, \dots, \alpha_n, \alpha_{n+1}), \dots, \alpha = (k_1, \dots, k_n, k_{n+1}).$$

Theorem VI. The sys. in (2) is solvable if and only if the vector α of V generated by the n vectors $\alpha_1, \dots, \alpha_n$ and the vector space X generated by the $n+1$ vectors $\alpha_1, \dots, \alpha_{n+1}$ have the same rank. If (2) is not solvable then $\text{rank } V = 1 + \text{rank } X$.

Proof. As every linear homogeneous relation between the α_i holds when l_1, \dots, l_n are replaced by 0, the rank r of V is not greater than the rank r' of X . Let X and V be the vector spaces generated by the solutions of (15) and (16) then $\text{rank } X = n - r'$, $\text{rank } V = n + 1 - r$. If (x_1, \dots, x_n) is a solution of (2) then $(x_1, \dots, x_n, 0)$ is a solution of (15). Therefore $\text{rank } X \leq \text{rank } V$, $r + 1 \leq n - r' \leq n - r$. If $r = n + 1$, $\text{rank } X = \text{rank } V$ then (15) has only solutions with $x_{n+1} = 0$ and (2) has r solutions. If $r = r'$, $\text{rank } X > \text{rank } V$ and there exists a solution

$(x_1, \dots, x_n, x_{n+1}) \neq 0$ of (15). Hence $\begin{pmatrix} x_1 \\ \vdots \\ x_n \\ x_{n+1} \end{pmatrix}$ is a solution of (2).

Theorem VII. If there are solutions of (2), we will get them by "sweep out".

Proof. We can get a basis of the solutions X of (15) by the method of sweep out. If there is a solution in which $x_{n+1} \neq 0$ it must also be a basis $(n+1)$ vector of this kind. From this $(n+1)$ vector we will get a solution of (2) on dividing the coordinates by x_{n+1} .

As 0 is a solution of (2) if and only if the equations are homogeneous the solutions form a vector space only in this case. In general the set of all solutions will be described by the following theorem.

Theorem VIII: If ξ and η are solutions of (2) and $s+t=1$ then $s\xi+t\eta$ is also a solution. Consequently if a set W of n vectors has the property that with two n -vectors ξ and η the n -vector $s\xi+t\eta$ belongs to W , then W is the set of the solutions of a system (2).

Proof: From $\sum_{i=1}^n a_{ij}x_j = a_i$, $\sum_{i=1}^n a_{ij}y_j = a_i$ it follows that $\sum_{i=1}^n a_{ij}(sx_j+ty_j) = a_i$. As the same holds for the other equations of \mathcal{A} the first proposition is proved. If χ, λ, π are n -vectors of W then $2(1/2\chi + 1/2\lambda - \pi) = \chi + \lambda - 2\pi$ and $\chi + \lambda - 2\pi$ belongs to W . Hence $(\chi - \pi) + (\lambda - \pi) = \chi + \lambda - 2\pi$ and $(\chi - \pi) + (\lambda - \pi)$ belongs to W . If π is a fixed element of W , the set of all differences $(\chi - \pi)$ is a vector space X of V . From Theorem III it follows that X is the set of the solutions of a system of linear homogeneous equations. This system may be (2.11). If $a = p_1, \dots, p_r$ and we define a_{ij} by $a_{ij} = \sum_{k=1}^r p_k x_{kj}$, $b_{ij} = \sum_{k=1}^r p_k y_{kj}$, then by Theorem I W is the set of all solutions of (2).

3. THE METHOD OF ORTHOGONALIZATION

By the previous theorems the problem proposed in the introduction has been solved completely. It will always be the principal part of the calculation to find out a basis of the vector space $X = W$. We may do it by the method of "sweep out" but it is important to have other suitable methods for it. In this section the method of orthogonalization will be treated; its advantage is that we will get at the same time a basis of X and a basis of V , both in a special form.

In the previous, and in the following paragraphs there is no restriction about the numbers which should be used.¹ However in this section we will suppose that all numbers used are real.

1) $\chi \neq 0$. The scalar product S of two n -vectors $\alpha = (\alpha_1, \dots, \alpha_n)$ and $\beta = (\beta_1, \dots, \beta_n)$ is:

$$S(\alpha, \beta) = S(\beta, \alpha) = \sum_{i=1}^n \alpha_i \beta_i. \quad (16)$$

Definition 10: The **Length** of α is a number $|\alpha| \geq 0$, defined by

$$|\alpha|^2 = S(\alpha, \alpha). \quad (17)$$

¹ (The same the restriction to real numbers introduced in Chapter II.) The theorems and the proofs hold for an arbitrary field of characteristic $\neq 2$.

Formulas $\alpha \cdot \beta = 0$ and on γ if $\alpha = 0$

$$\mathcal{S}(\alpha + \gamma) = \mathcal{S}\alpha + \mathcal{S}\gamma$$

$$\mathcal{S}\alpha\beta = \alpha\mathcal{S}\beta$$

(18)

$$|\alpha + \beta|^2 \leq |\alpha|^2 + |\beta|^2 \quad \text{Cauchy's inequality,}$$

$$|\alpha + \beta| \leq |\alpha| + |\beta|$$

Definition 11. If $\mathcal{S}\alpha = 0$, α and β are orthogonal

Lemma 1. In the next section we shall explain some of these notions and formulas will be explained

The coordinates of α can be expressed as scalar products

$$\alpha_k = \mathcal{S}\alpha e^k \quad \dots (19)$$

Definition 12. The vectors β^1, \dots, β^n form an orthogonal system if they satisfy the conditions:

$$\mathcal{S}\beta^i \beta^i = 1$$

$$\mathcal{S}\beta^i \beta^j = 0 \quad i \neq j. \quad \dots (20)$$

Properties of orthogonal systems

1. If the β^i form an orthogonal system they are independent

Proof. From $0 = \sum c_i \beta^i$ it follows

$$0 = \beta^k \sum c_i \beta^i = c_k \text{ for } k=1, \dots, n.$$

2. If the β^i form an orthogonal system, and α is independent of the β^i , then there is an n vector β^{n+1} such that α is dependent on $\beta^1, \dots, \beta^{n+1}$, and these $n+1$ vectors form an orthogonal system.

Proof. Let $\alpha = \sum \mathcal{S}\alpha \beta^i \beta^i$, and $\beta^{n+1} = |\alpha|^{-1} \alpha$, then there is for $k=1, \dots, n$, $0 = \mathcal{S}\beta^k \alpha = \mathcal{S}\beta^k \sum \mathcal{S}\alpha \beta^i \beta^i = \mathcal{S}\beta^k \beta^i \mathcal{S}\alpha = 0$

3. If V is a vector space of rank m containing a vector space A of rank $r < m$ then there exists a basis of V forming an orthogonal system β^1, \dots, β^m , such that β^1, \dots, β^r form a basis of A .

Proof. Let β^1 be an arbitrary n vector of V of the length 1. If $r > 1$ there is in A an n vector independent of β^1 , and therefore there is an orthogonal system β^1, β^2 . By repeated use of this construction we get β^1, \dots, β^r in A and by continuing the method in V we will get $\beta^{r+1}, \dots, \beta^m$.

Theorem IX. Let V be the vector space generated by the vectors of the equations (2.11) and X the vector space of the solutions of (2.11). Then there exists an orthogonal system $a^1, \dots, a^r, \xi^1, \dots, \xi^{n-r}$ such that a^1, \dots, a^r form a basis of V , and ξ^1, \dots, ξ^{n-r} form a basis of X .

Proof. The connection between the vector spaces V and X is that every vector of X is orthogonal to every vector of V and $r > 0$ and $n > 0$. Hence if we construct an orthogonal system of n vectors, $a^1, \dots, a^r, \xi^1, \dots, \xi^{n-r}$, such that the first r vectors form a basis of V , the ξ^1, \dots, ξ^{n-r} are orthogonal to the r vectors of V and form therefore solutions of (2.11). An arbitrary vector can be expressed by $\lambda = \sum_{i=1}^r a^i + \sum_{j=1}^{n-r} \xi^j$. λ is orthogonal to the r vectors of V and only if for $k=1, \dots, r$ $0 = \sum_{i=1}^r \lambda_{i,k} = \sum_{j=1}^{n-r} b_{j,k} \lambda_{j,k}$. ξ^1, \dots, ξ^{n-r} is therefore a basis of X .

Remark. The proof of Theorem IX does not apply to the Theorems IV-VIII and the method of sweep out. Theorem IV follows directly from Theorem IX.

§ 10. SUBSTITUTION AND ELIMINATION

The method of Substitution. In order to get the solutions of (2) we can also reduce the problem to $m-1$ equations and $n-1$ unknown. Let x_{i_1} be the first unknown for which the coefficient of the 1st equation does not vanish then:

$$x_{i_1} = - \sum_{j=1}^n a_{1,j} x_{i_1,j} + \frac{a_{1,i_1}}{a_{1,i_1}}$$

On substituting this value in the other $m-1$ equations we get $m-1$ equations with $n-1$ unknown. If some of these equations become identical they will be omitted. The rows of the new equations are

$$(l_i = (i_1, k) : a_{i,k}(1).$$

Hence the substitution of x_{i_1} is nothing more than the sweep out of the column $\langle i_1 \rangle$. On continuing this procedure the column $\langle i_2 \rangle$ will be swept out in the rows l_1, \dots, m also $\langle i_1 \rangle$ in the rows $(3), \dots, m$ etc. the 0 rows, as they belong to identical equations being always omitted. The result will be that the columns $\langle i_1 \rangle, \dots, \langle i_r \rangle$ will be swept out in the rows below $(i_r + 1)$. The last of these equations give us the possibility to express x_{i_r} by the remaining $n-r$ unknown. The substitution of this

value in the i -th position is nothing but the sweep-out of the column $\langle i \rangle$ in the rows above the row i . On repeating this procedure by the substitution of x_{i+1}, \dots, x_n expressed by the last $n-i$ unknown we sweep out the corresponding column in the upper rows. Hence the method of substitution is not quite different from the method of sweep-out.

Theorem 1.1. Proposition. The equations (2.11) will be solved when we have a partition of the set $\{1, \dots, n\}$ into vectors belonging to (1.1) are independent on the vectors belonging to (2.11). We should multiply the equations (2.11) with suitable numbers and add so that the sum of these equations becomes an equation of the kind (1.1). We eliminate the unknown x_1, \dots, x_r to get the first row of (1.1), x_{r+1}, \dots, x_n to get the second row, etc. By this procedure we intend to find out at once what is the step by step on using the method of sweep-out. The rank r and also the factors we need for the elimination are dependent only on the given coefficients. We have to find out a character of that dependence. For this purpose we need the determinants.

11. DEFINITION AND PROPERTIES OF THE DETERMINANTS

Definition 1.1. A function

$$D = D(a^1, \dots, a^n) \quad (2.1)$$

of n vectors

$$a^1, \dots, a^n$$

is called a *Determinant*

$$D = D(a^1, \dots, a^n) \quad (2.2)$$

$$= \det a_i^j$$

if the following conditions are satisfied:

$$(a) \quad D(a^1, \dots, a^i, \dots, a^i, \dots, a^n) = 0, \quad \text{if } i \neq j \text{ for } i = 1, \dots, n$$

$$(b) \quad D \text{ will not change by replacing}$$

$$a_i \longrightarrow a_i + a_j, \text{ for } i \neq j, \quad \dots \quad (2.3)$$

$$(c) \quad D(a^1, \dots, a^n) = 1.$$

It will be proved later on that a function of this kind exists and that it is uniquely defined by the conditions (a), (b), (c). Now we will consider

the properties of a function $n+1$ to be satisfied simultaneously the above conditions

1. If $a^i = 0$, $L = 0$.

Proof. $0 = 0 \cdot L = L(a^1, \dots, 0, a^i, \dots, a^n) = L, a^i = 0$.

2. L will not be changed by replacing $a^i \rightarrow a^i + c$ if $c^i \in F_{a^i}$ and c is an arbitrary number

Proof. $L = \frac{1}{n!} \frac{\partial^n L}{\partial a^1 \dots \partial a^n} (a^1, \dots, a^n)$
 $= \frac{1}{n!} \frac{\partial^n L}{\partial a^1 \dots \partial a^n} (a^1, \dots, a^i + c, \dots, a^n)$
 $= \frac{1}{n!} \frac{\partial^n L}{\partial a^1 \dots \partial a^n} (a^1, \dots, a^i, \dots, a^n)$

3. On differentiating L at $a^i = 1$ we get equality 1.

Proof. $L = L(a^1, \dots, a^i, \dots, a^n)$

$= \frac{1}{n!} \frac{\partial^n L}{\partial a^1 \dots \partial a^n} (a^1, \dots, a^i, \dots, a^n)$
 $= \frac{1}{n!} \frac{\partial^n L}{\partial a^1 \dots \partial a^n} (a^1, \dots, a^i, \dots, a^n)$
 $= \frac{1}{n!} \frac{\partial^n L}{\partial a^1 \dots \partial a^n} (a^1, \dots, a^i, \dots, a^n)$

4. If a^i is dependent on the other $n-1$ then $a^i = 0$ and $L = 0$

Proof. From condition 3 that a^i can be replaced by 0. Hence $L = 0$.

$L(a^1, \dots, a^n) = 1$ if a^1, \dots, a^n is an even permutation
 $= -1$, if a^1, \dots, a^n is an odd permutation,
 $= 0$, if the indices are not different.

These formulae follow directly from (1), (3) and (4)

B. $L(a^1, \dots, a^i + y, \dots, a^n) = L(a^1, \dots, a^n) + y \frac{\partial L}{\partial a^i}$

Proof. Either a^i is dependent on the other n vectors a^1, \dots, a^n or these $n+1$ vectors are not independent or the n vectors a^1, \dots, a^n are all independent. In the first case, $L=0$ on the left side a^i may be omitted and therefore the equation holds. In the 2nd case each of the three functions equals zero. In the 3rd case a^i is dependent on the other $n-1$ vectors. By applying 2 on both sides of the equation the left side becomes $(1+y) L$, and the right side becomes $L + y \frac{\partial L}{\partial a^i}$, therefore it holds.

7. If

$$L_1 = \sum_{j=1}^n a_{1j} x_j, \quad \text{and } D_1 \text{ is the 1st determinant we get by}$$

replacing $x_1 \rightarrow \sum_{j=1}^n a_{1j} x_j$ (being not clear as to $i=1$ or i then $L_1 = \sum_{j=1}^n a_{1j} x_j$),

Proof. To prove the theorem we have to use 6 and 7. Since

8. Let L_1 be the determinant we get by replacing x_1 by x_1 and

$$\text{let } a^1 = \sum_{j=1}^n a_{1j} x_j, \text{ then}$$

$$L_1 = \sum_{j=1}^n a_{1j} L_j \quad \dots \quad (24)$$

The theorem follows directly from 7.

$$9. L = \sum_{j=1}^n a_{1j} x_j \dots x_n \quad \dots \quad (25)$$

the summation has to be made only for the permutations $\sigma_1, \dots, \sigma_n$ if n is odd and for all permutations if n is even.

$$L_1 = \sum_{j=1}^n a_{1j} x_j \dots x_n \quad \dots \quad (26)$$

hence from 8 follows the theorem.

Theorem X. The determinant $D(x_1, \dots, x_n)$ exists and is uniquely defined by every ordered set of n vectors a^1, \dots, a^n .

Proof. As it was stated in 9 the value of the determinant—if there exists one—cannot be different from 1. Therefore we have to prove that (25) has the properties (a), (b), (c), (d), and (e) are consequently satisfied we will prove (d). The right side of (25) is a function of x_1, \dots, x_n (may be called $D(x^1, \dots, x^n)$). By interchanging two of the vectors of the argument the even permutations become odd and vice versa, therefore D is changed to $-D$. Hence if $x_1 = x_2$ then $D(x_1, x_2, \dots, x_n) = 0$. From (25) it follows that $D(x^1, \dots, x^n) = \sum_{j=1}^n a_{1j} x_j \dots x_n$ and as the last vector vanishes (b) is satisfied and the theorem holds.

§ 12. Further Properties of the Determinant

The columns of (22) form a vector

$$a_1 = (a_{11}, a_{21}, \dots, a_{n1})^T \quad (27)$$

$$a_2 = (a_{12}, a_{22}, \dots, a_{n2})^T$$

$$I = I(a_1, \dots, a_n) = I(a_1, \dots, a_n) \quad (23)$$

i.e., a determinant which is unchanged if the rows are interchanged by the transposition.

Let $I(a_1, \dots, a_n)$ be the polynomial inverse to $I(a_1, \dots, a_n)$, then these polynomials are either both even or both odd.

$$I(a_1, \dots, a_n) = I(a_1, \dots, a_n) \quad \text{and} \\ I(a_1, \dots, a_n) = I(a_1, \dots, a_n) \quad (24)$$

Let $I(a_1, \dots, a_n)$ be the determinant $I(a_1, \dots, a_n)$ if we replace the a_i by the a_i .

Then the above are even or odd functions of (a_i) .

11. If we replace a_i by a_i the determinant becomes independent of the co-ordinates of a_i .

Proof. On replacing a_i by a_i the product $a_1 \dots a_n$ becomes $a_1 \dots a_n$ hence a certain number of terms in the determinant are replaced by 1 hence the determinant becomes independent of the co-ordinates of a_i .

If $I(a_1, \dots, a_n)$ we get on replacing a_i by a_i , $a_i = p_1, \dots, p_n$ being different and $a_i = q_1, \dots, q_n$ being different will be called

$$I(p_1, \dots, p_n) \quad (25)$$

12. Let $I(a_1, \dots, a_n)$ be the determinant we get on replacing in $I(a_1, \dots, a_n)$ the arguments a_1, \dots, a_n by a_1, \dots, a_n then

$$I = I(a_1, \dots, a_n) \quad \text{hence}$$

Proof. From (10) it follows that I is not changed if we replace the a_i of (10) by a_i from (10) and (11) it follows also that I is not changed if we replace the a_i of (10) by a_i . In both cases we get a determinant with

π_1, \dots, π_n is an arbitrary permutation of $1, \dots, n$, and τ_1, \dots, τ_n an arbitrary permutation of $1, \dots, n$. The sign \pm has to be taken if and only if both permutations are of the same kind \pm or $-$, if

$$u_1, \dots, u_m, v_1, \dots, v_{n-m} \quad (33)$$

is an even permutation of $1, \dots, n$,

and therefore also an even permutation of $\tau_1, \dots, \tau_m, \tau_1, \dots, \tau_{n-m}$.

Every permutation of $1, \dots, n$ has $n!$ the number of the monomials in Δ is $n!$ and therefore the sum becomes L . In the case of an odd permutation the sign of every monomial \pm is changed and therefore we get $-L$. When the numbers j and r are not all different Δ becomes a determinant with two identical rows and from (10.4) it follows, that it is vanishing.

The relation between the two factors of each term of Δ is a reciprocal one: they are called *cofactors*. We can extend the notion of cofactor and the formula (33) also to the case $n=m$ by taking

$$L_1^1 = a_1^1, L_2^2 = a_2^2, \dots, L_n^n = a_n^n, \quad L_1^r = -a_1^r, \dots, L_n^r = -a_n^r$$

Often it is used to denote $m=n-1$ and

$$\sum_{j=1}^n a_j^i L_j^i = 0, \text{ for } i \neq j \\ = L_i, \text{ for } i=j. \quad (34)$$

From (34) it follows that (34) holds also when the upper and the lower indices are not exchanged: hence we get

$$\sum_{j=1}^n a_j^i L_j^i = 0, \text{ for } i \neq j \\ = L_i, \text{ for } i=j. \quad (34')$$

If an ordinary Matrix D has the elements $[a_{ij}] = a_{ij}$ we will write

$$D = (a_{ij})$$

On selecting some rows and some columns of D , we get new matrices

$$D_{i_1, \dots, i_r}^{j_1, \dots, j_r} = (a_{i_k j_k}), \quad a_{i_k}^{j_k} = a_{i_k}^{j_k}$$

If $n = 1$, $m = n$ the matrix is defined to be a unit, and the Minor of D of order m

$$\det D_{h_1, \dots, h_m}^{h_1, \dots, h_m} = \det \delta_{h_i}^{h_i} \quad (11a)$$

13. If $L_{h_1, \dots, h_m}^{h_1, \dots, h_m} P_r$ and $L_{h_1, \dots, h_m}^{h_1, \dots, h_m}$ are cofactors then

$$L_{h_1, \dots, h_m}^{h_1, \dots, h_m} P_r = \det D_{h_1, \dots, h_m}^{h_1, \dots, h_m}$$

Proof. The replacement P_r by which we get $L_{h_1, \dots, h_m}^{h_1, \dots, h_m}$ from L , annihilate in L every term not having the factor $\delta_{h_i}^{h_i}$ and in the remaining monomials the factor will be replaced by 1. The monomials of $L_{h_1, \dots, h_m}^{h_1, \dots, h_m}$ are therefore just the monomials of a determinant composed of the rows r_1, \dots, r_{n-m} and the columns $\langle i_1 \rangle, \dots, \langle i_{n-m} \rangle$. For getting the proper sign we must realize that the permutation of the $n-m$ rows is of the same character even or odd as the permutation of the $n-m$ columns, if and only if the permutation of the original n rows is of the same character as the permutation of the original n columns. This condition is however satisfied by the definition of the cofactors.

14. If all minors of D of order m vanish then the minors of higher order vanish also.

Proof. Using the formula (11) we can develop an arbitrary minor of order $m+1$ as a homogeneous linear function of degree m with $n-m$ coefficients being minors of order m .

Definition 14. If there is a non vanishing minor of D of order r but every minor of order $r+1$ vanishes r is called the rank of D .

17. The matrix D the vector space generated by its rows and the vector space generated by its columns have the same rank.

Proof. If r is the rank of the vector space of the row vectors there is a non homogeneous relation between every set of $r+1$ row vectors and this relation holds also if we get away some of the columns. Hence from n to $n-r$ columns, that the minors of order $r+1$ are all vanishing out

the rank of D is at most r . If we let d be the rank of D every minor of order $> d$ vanishes, and hence r of order r or R is not vanishing. Without restriction of the generality we suppose $D = (d_{ij})_{1 \leq i \leq R, 1 \leq j \leq n}$. Let (α_i) be an arbitrary row

In $D = (d_{ij})_{1 \leq i \leq R, 1 \leq j \leq n}$ the cofactors of $d_{i+1,1}, \dots, d_{i+1,n}$ will be

$$\text{called respectively } A_1, \dots, A_n, \quad A_i = D_{i+1,1}^{(i)} \dots d_{i+1,n}^{(i)} \neq 0$$

Hence the coordinates of $\alpha = \alpha^1 A_1 + \dots + \alpha^n A_n + \alpha^{n+1} A_{n+1}$ are either determinants with two equal rows or minors of D of order $R+1$ and therefore vanishing. Hence $\alpha = 0$, i.e., every α^i depends on $\alpha^1, \dots, \alpha^n$, hence the rank of the vector space is at most R . As we have seen in the first part of the proof the rank of the vector space is not smaller than R and therefore both ranks are equal. The rank of D is not changed, if the rows and the columns of D are interchanged; hence the rank of D equals n to the rank of the vector space generated by the columns of D .

* 18. PROPOSITION ON DETERMINANTS

Theorem XI. The equations (2) are equivalent and only if the matrices

$$M = \begin{pmatrix} a_1 & a_2 \\ b_1 & b_2 \end{pmatrix} \quad \text{and} \quad \tilde{M} = \begin{pmatrix} a_1 & a_2 & a_3 \\ b_1 & b_2 & b_3 \end{pmatrix} \quad (3)$$

have the same rank

Proof. From (3) it follows that the ranks of the matrices are the same as the ranks of the vector spaces V and \tilde{V} of Theorem VI. Hence the Theorem XI follows from the Theorem VI.

In order to get the rank r of a matrix by the help of determinants it is not necessary to calculate each minor, it is sufficient to state that one minor of order r is not vanishing and that all minors of order $r+1$ vanish. The rank does not change by row addition or by column addition, it is often useful to simplify matrices by these operations.

Corollary. If the matrices M and \tilde{M} of (3), have the same rank then the matrix formed by an arbitrary set of rows i_1, \dots, i_p of M has the same rank as a matrix formed by the same rows of \tilde{M} .

Proof: As the non-homogeneous equations belonging to M and \bar{M} are solvable, the rows $(p_1), \dots, (p_r)$ define solvable equations, and therefore the two matrices defined by these r rows have the same rank.

Theorem XII. If the matrices M and \bar{M} of Theorem XI have the same rank r , and a non-vanishing minor of \bar{M} has the rows $(p_1), \dots, (p_r)$, then the solutions of (2) are identical with the solutions of the equations belonging to these rows.

Proof: From the Corollary it follows that the rank of the matrix formed by the $(p_1), \dots, (p_r)$ is also r . The solutions of the homogeneous equations defined by the r rows form a vector-space X of rank $n - r$, including all solutions of (2/II). As the rank of the vector-space of the solutions of (2/II) is also $n - r$, it follows from §4, 7 that every solution of the r linear homogeneous equations is also a solution of (2/II). Hence from Theorem I it follows that we get every solution of (2) by adding all vectors of X to an arbitrary solution of the r non-homogeneous equations, i.e., the solutions of these equations are just the solutions of (2).

By the Theorems XI and XII the problem of solving a system of linear equations by elimination has been reduced to the following: Find out the solutions of

$$\begin{aligned} a_1^1 x_1 + \dots + a_r^1 x_r + a_{r+1}^1 x_{r+1} + \dots + a_n^1 x_n &= a_0^1 \\ a_1^2 x_1 + \dots + a_r^2 x_r + a_{r+1}^2 x_{r+1} + \dots + a_n^2 x_n &= a_0^2 \end{aligned} \quad (37)$$

when the determinant $\det A$ of the matrix

$$A = \begin{pmatrix} a_1^1, \dots, a_r^1 \\ \vdots \\ a_1^r, \dots, a_r^r \end{pmatrix} \quad \text{is not vanishing.}$$

Theorem XIII. Let A_k^i be the co-factors of a_k^i in A , then we get the solutions of (37) by

$$\det A \cdot x_k + \sum_{i=1}^r a_{i+1}^k A_i^k x_{i+1} + \dots + \sum_{i=1}^r a_n^k A_i^k x_n = \sum_{i=1}^r a_0^i A_i^k \quad (38)$$

$$k=1, \dots, r.$$

Proof: By multiplying the equations (37) respectively by the co-factors A_k^i (k being constant), and adding, we get the equations (38).

These equations are therefore necessary conditions for the solutions of (37). The rows of (38)—including the right sides of the equations—are dependent on the rows of (37). The rank of the matrix composed of the rows of (37) and (38) is therefore also r . The determinant formed by the first r columns of (38) is a power of $\det A$ and therefore $\neq 0$. Hence from Theorem XII it follows that the solutions of (38) are identical with the system of all solutions of (37) and (38), and therefore (38) is also a sufficient condition for the solutions of (37).

§ 14. LINEAR TRANSFORMATIONS.

Let

$$A = \begin{pmatrix} a_1^1 & \dots & a_n^1 \\ \dots & \dots & \dots \\ a_1^n & \dots & a_n^n \end{pmatrix} \quad (39)$$

be a matrix with n rows and n columns.

The row-vectors are called $\alpha^1, \dots, \alpha^n$,

the column-vectors are called $\alpha_1, \dots, \alpha_n$.

We consider the equations

$$\sum_{j=1}^n a_i^j x_j = y_i, \quad i=1, \dots, n. \quad (40)$$

To every n -vector $\xi = (x_1, \dots, x_n)$ corresponds an n -vector $\eta = (y_1, \dots, y_n)$;

(40) is called a *linear transformation*, and we will express it by

$$\xi \longrightarrow \eta.$$

Theorem XIV. By a linear transformation (40) the n -vectors ξ are transformed to the n -vectors of a vector-space, whose rank equals the rank of A .

Proof: If $\xi^1 \longrightarrow \eta^1$, $\xi^2 \longrightarrow \eta^2$, then $\xi^1 + \xi^2 \longrightarrow \eta^1 + \eta^2$

and $c\xi^1 \longrightarrow c\eta^1$, for every number c .

Hence from §4, 10 it follows, that η form a vector-space H . Two n -vectors ξ are transformed to the same n -vector η , if and only if the difference is transformed to 0, i.e., if the difference belongs to the vector-space Z of the solutions of the homogeneous system belonging to (40). Using the methods of §4, 7, or §8, 3, we will find out a basis ξ^1, \dots, ξ^{n-r} , ξ^1, \dots, ξ^r of the vector-space of ξ , so that the $n-r$ first n -vectors form

a basis of Z . If $\xi^i \rightarrow \eta^i$, an arbitrary n -vector is transformed $\sum c_i \xi^i + \sum d_{ij} \zeta^j \rightarrow \sum c_i \eta^i$, and from the definition of ζ^i it follows that this n -vector vanishes if and only if $c_1 = \dots = c_r = 0$. Therefore η^1, \dots, η^r form a basis of H , and r is the rank of H .

Theorem XV. A representation of the n -vectors ξ by the n -vectors $f(\xi)$ with the properties: $f(\xi^1 + \xi^2) = f(\xi^1) + f(\xi^2)$, $f(c\xi) = cf(\xi)$, is a linear transformation.

Proof. Let $f(\xi^i) = \beta_i = \sum b_{ij} \zeta^j$, $\xi = \sum x_i \xi^i$, then $f(\xi) = \sum b_{ij} x_i \zeta^j$.

Hence the representation is a linear transformation with the matrix $((b_{ij}))$.

If $y_i = \sum a_{ij} x_j$, $x_i = \sum b_{ij} z_j$, then $y_i = \sum a_{ij} b_{jk} z_k = \sum d_{ik} z_k$, where

$$d_{ik} = \sum a_{ij} b_{jk} = B \cdot A' \beta_i, \quad (41)$$

α' being the row-vectors of $A = ((a_{ij}))$, and β_j the column-vectors of $B = ((b_{ij}))$. The matrix $D = ((d_{ik}))$ is called the product

$$D = A \cdot B \quad (42)$$

If we consider 3 matrices A, B, C , and their products $(A \cdot B) \cdot C$ and $A \cdot (B \cdot C)$, we get in both cases a matrix with the co-ordinates $d_{ik} = \sum_{j=1}^n \sum_{l=1}^n a_{ij} b_{jl} c_{lk}$, hence:

Theorem XVI. The associative law holds for the multiplication of matrices.

§15. DECOMPOSITION OF MATRICES.

A matrix $D = ((d_{ij}))$, $d_{ii} = d_i$, $d_{ij} = 0$, when $i \neq j$, is called a *Diagonal-matrix*. A matrix $E_{rs}(\lambda) = ((e_{ij}))$, $e_{ii} = 1$, $e_{rs} = \lambda$, and every other $e_{ij} = 0$, is called an *Elementary-matrix* (defined only for $r \neq s$). If A has the row-vectors $\alpha^1, \dots, \alpha^n$, and the column-vectors $\alpha_1, \dots, \alpha_n$, then

$D \cdot A$ has the row-vectors $d_1 \alpha^1, \dots, d_n \alpha^n$

$A \cdot D$ has the column-vectors $d_1 \alpha_1, \dots, d_n \alpha_n$

$E_{rs}(\lambda) \cdot A$ has the row-vectors $\beta^i = \alpha^i$,

$$\text{for } i \neq r \quad \beta^r = \alpha^r + \lambda \alpha^s \quad (48)$$

$A \cdot E_{rs}(\lambda)$ has the column-vectors $\gamma_i = \alpha_i$,

$$\text{for } i \neq s \quad \gamma_s = \alpha_s + \lambda \alpha_r \quad \checkmark$$

Theorem XVII. An arbitrary matrix A of n rows and n columns is a product $A = P_1 \cdot D \cdot P_2$, where P_1 and P_2 are products of elementary-matrices, and D is a diagonal-matrix.

Proof. From (43) it follows that the theorem is identical with the following: A can be transformed to a diagonal-matrix by row-additions $\alpha' \rightarrow \alpha' + \lambda \alpha''$, and by column-additions $\alpha' \rightarrow \alpha' + \lambda \alpha''$. In order to prove this proposition, we use the method of "sweep-out" in a little modified manner. If every element of A vanishes, A is a diagonal-matrix; if all elements do not vanish, we can make $[1, 1] \neq 0$ by row-additions and column-additions of the type mentioned above, and by the same kind of operations we can sweep out the first row and the first column. On continuing this procedure we get a diagonal-matrix.

Theorem XVIII. The determinant of $A \cdot B$ is equal to the product of the determinants of A and B , i.e.,

$$\det(A \cdot B) = \det A \cdot \det B.$$

Proof. From (43) it follows, that the determinant of a matrix does not change, when the matrix is multiplied with an elementary-matrix, and therefore the determinant does not change, when the matrix is multiplied with a product of elementary-matrices. Hence, if $B = P_1 \cdot D \cdot P_2$, $\det B = \det D = d_1 \dots d_n$. On the other hand $\det(A \cdot B) = \det(A \cdot P_1 \cdot D)$ holds. As we get $\frac{1}{d_i} P_1 \cdot D$ from $A \cdot P_1$ by multiplying the columns with $d_1 \dots d_n$, we get: $\det(A \cdot B) = \det(A \cdot P_1 \cdot D) = \det(A \cdot P_1) \cdot d_1 \dots d_n = \det A \cdot \det B$.